The Ostrovsky-Vakhnenko equation on the half-line: a Riemann-Hilbert approach

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Abstract

We analyze an initial-boundary value problem for the Ostrovsky-Vakhnenko equation

$$u_{xxt} - 3u_x + 3u_x u_{xx} + u u_{xxx} = 0$$

on the half-line. This equation can be viewed as the short wave model for the Degasperis-Procesi (DP) equation. We show that the solution u(x,t) can be recovered from its initial and boundary values via the solution of a 3×3 vector Riemann-Hilbert problem formulated in the complex plane of a spectral parameter z.

Key Words: Ostrovsky-Vakhnenko equation, initial-boundary value problem, Riemann-Hilbert problem

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1 Introduction

In this paper, we consider the partial differential equation (PDEs)

$$u_{xxt} - 3bu_x + 3u_x u_{xx} + uu_{xxx} = 0 (1.1)$$

where b > 0 is a parameter and u = u(x,t) is real-valued. This equation stems from the short-wave limit of the Degasperis-Procesi equation (DP) [1], which is a model of nonlinear shallow water waves:

$$u_t - u_{xxt} + 3\omega u_x + 4uu_x = 3u_x u_{xx} + uu_{xxx}. (1.2)$$

Indeed, introducing new space-time variables (x', t') and a scaling of u by

$$x' = \frac{x}{\varepsilon}, \quad t' = \varepsilon t, \quad u' = \frac{u}{\varepsilon^2},$$

where ε is a small positive parameter, then (1.1) is the leading term of (1.2) as $\varepsilon \to 0$. Thus, equation (1.1) can be named as the short wave model of the Degasperis-Procesi equation.

Equation (1.1) arises also in the theory of propagation of surface waves in deep water, see [2], as an asymptotic model for small-aspect-ratio waves.

For $\omega = 0$, equation (1.1) reduces to the derivative Burgers equation

$$(u_t + uu_r)_{rr} = 0,$$

whereas for $b=-\frac{1}{3}$, it reduces to the (differentiated) Vakhnenko equation [3, 4]

$$(u_t + uu_x)_x + u = 0. (1.3)$$

Alternatively, (1.1) with $b = \frac{1}{3}$ reduces to (1.3) after the change of variables $(u,t) \to (-u,-t)$.

Equation (1.3) arises—and is known as the Vakhnenko equation in the context of propagation of high-frequency waves in a relaxing medium [3, 5, 6]. On the other hand, being written in the form

$$(u_t + c_0 u_x + \alpha u u_x)_x = \gamma u, \tag{1.4}$$

it is also called the reduced Ostrovsky equation [7]: it corresponds, in the case $\beta = 0$, to the equation

$$(u_t + c_0 u_x + \alpha u u_x + \beta u_{xxx})_x = \gamma u, \tag{1.5}$$

that was derived by Ostrovsky in 1978 [8]. Therefore, it is more correct to name equation (1.3) the OstrovskyCVakhnenko equation (OV), as it is proposed in [9].

Well-posedness of the Cauchy problem for the Ostrovsky equation and its relatives (reduced Ostrovsky equation, generalized Ostrovsky equation, etc) in Sobolev spaces has been widely studied in the literature, using PDE techniques; see [10, 11, 12, 13, 14].

On the other hand, equation (1.1) is (at least, formally) integrable: it possesses a Lax pair representation (see [15] appendix, (A1))

$$\psi_{xxx} = \lambda(-u_{xx} + b)\psi \tag{1.6a}$$

$$\psi_t = \frac{1}{\lambda}\psi_{xx} - u\psi_x + u_x\psi \tag{1.6b}$$

where $\psi = \psi(x, t, \lambda)$.

It is because of integrable, the Ostrovsky equation can be solved by the inverse scattering method. Recently, Boutet de Monvel and Shepelsky formulate a Riemann-Hilbert problem to the initial value problem, and they also consider the long-time asymptotic problem, see [16].

However, in many laboratory and field situations, the wave motion is initiated by what corresponds to the imposition of boundary conditions rather than initial conditions. This naturally leads to the formulation of an initial-boundary value (IBV) problem instead of a pure initial value problem.

In 1997, Fokas announced a new unified approach for the analysis of initial-boundary value problems for linear and nonlinear integrable PDEs [17, 18, 19]. The unified method provides a generalization of the inverse scattering formalism from initial value to initial-boundary value problems.

Over the last almost two decades, it has been used to analyze boundary value problems for several important integrable equations with 2×2 Lax pairs. Recently, Lenells develop a methodology for analyzing initial-boundary value problems for integrable evolution equations with Lax pairs involving 3×3 matrices [20]. Then many other integrable evolution equations with 3×3 Lax pairs are analyzed too, see [21, 22, 23].

In this paper, we use Fokas and Lenells method to analyze the initial-boundary value problem for the Ostrovsky-Vakhnenko equation (1.1) on the half-line domain, that is, in the domain

$$\Omega = \{(x, t) \in \mathbb{R}^2 | 0 \le x < \infty, 0 \le t < T\},\tag{1.7}$$

where $T < \infty$ is a given positive constant. Assuming that a solution exists, we show that u(x,t) can be recovered from the initial and boundary values $u_0(x)$, $g_0(t)$, $g_1(t)$, $g_2(t)$ defined by

$$u(x,0) = u_0(x), \quad 0 < x < \infty,$$

 $u(0,t) = g_0(t), \quad u_x(0,t) = g_1(t), \quad u_{xx}(0,t) = g_2(t), \quad 0 < t < T.$ (1.8)

The main peculiarity compared with other applications of the approach of [17] is that the Lax pair involves 3×3 matrices instead of 2×2 matrices. This difference leads to some new challenges. Apart from the 3×3 Lax pair, the spectral analysis of equation (1.1) on the half-line also presents some other peculiarities: (a) The presence of singularities in the Lax pair implies that it is necessary to introduce two sets of eigenfunctions. The eigenfunctions in the first set are well-behaved near z = . The eigenfunctions in the second set are well-behaved near $z = \infty$. Together these two sets of eigenfunctions can be used to formulate a Riemann-Hilbert problem. An analogous situation occurs in the analysis of DP equation on the half-line in [21]. (b) The basic matrix eigenfunctions which are natural candidates for the formulation of a Riemann-Hilbert problem, are difficult to recover

the solution of the PDEs (2.1). We overcome this problem by formulating an associated vector Riemann-Hilbert problem, for which it is much more easier to recover the solution of our problem in hand. (c) The formulation of the Riemann-Hilber problem depends, in addition to the variables (x,t), on a function y(x,t) which is unknown from the point of view of the inverse problem. In order to obtain a Riemann-Hilbert problem whose jump matrix involves only known quantities, we have to reparametrize the x variable. This implies that we only obtain a parametric representation for the solution u(x,t).

We will consider the initial-boundary value problems for (1.1) for which the initial and boundary values satisfy

$$-u_{0xx}(x) + b > 0, \quad x \ge 0,$$

-g₂(t) + b > 0, $0 \le t < T$, (1.9)

as well as

$$g_0(t) \le 0, \quad 0 \le t < T.$$
 (1.10)

The assumptions in (1.9) imply that

$$-u_{xx}(x,t) + b > 0, \quad 0 \le x < \infty, 0 \le t < T.$$
 (1.11)

The assumption (1.10) is used to ensure boundedness of certain eigenfunctions.

The organization of the paper is as follows. In the following section 2, we introduce two sets eigenfunctions. In section 3, we derive expressions for the jump matrices in terms of suitable spectral functions. In section 4, we derive residue conditions for the pole singularities of the eigenfunctions. In section 5, we state our main result, see Theorem 5.1.

2 Spectral Analysis

In this section, starting from the Lax pair of Ostrovsky-Vakhnenko equation, see (1.6), we define analytic eigenfunctions which are suitable for the formulation of a Riemann-Hilbert problem.

Without loss of generality, in what follows we assume that b = 1. That is, we consider the the initial and boundary problems of the following equation

$$u_{xxt} - 3u_x + 3u_x u_{xx} + uu_{xxx} = 0 (2.1)$$

2.1 Two sets of eigenfunctions

Let z be the spectral parameter defined by $\lambda=z^3$. The coefficients of the original Lax pair (1.6) have singularities at $z=\infty$ and also at z=0. In order to have a good control on the behavior of eigenfunctions at $z=\infty$ and at z=0 we introduce new forms of (1.6), the first one appropriate at $z=\infty$, the second one at z=0.

Denote

$$\Lambda(z) = \begin{pmatrix} \lambda_1(z) & 0 & 0 \\ 0 & \lambda_2(z) & 0 \\ 0 & 0 & \lambda_3(z) \end{pmatrix},$$
(2.2)

where $\lambda_j(z) = z\omega^j$, j = 1, 2, 3, here $\omega = e^{\frac{2\pi i}{3}}$.

2.1.1 Lax pair (well-controlled at $z = \infty$)

Let
$$\Psi = \begin{pmatrix} \psi \\ \psi_x \\ \psi_{xx} \end{pmatrix}$$
, then the Lax pair (1.6) can be written as

$$\begin{cases}
\Psi_{x} = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ z^{3}q^{3} & 0 & 0 \end{pmatrix} \Psi, \\
\Psi_{t} = \begin{pmatrix} u_{x} & -u & z^{-3} \\ 1 & 0 & -u \\ -z^{3}uq^{3} & 1 - u_{x} \end{pmatrix} \Psi,
\end{cases} (2.3)$$

where $q^{3}(x,t) = -u_{xx}(x,t) + 1$.

Denote

$$D(x,t) = \begin{pmatrix} q^{-1}(x,t) & 0 & 0\\ 0 & 1 & 0\\ 0 & 0 & q(x,t) \end{pmatrix},$$
 (2.4a)

$$P(z) = \begin{pmatrix} 1 & 1 & 1\\ \lambda_1(z) & \lambda_2(z) & \lambda_3(z)\\ \lambda_1^2(z) & \lambda_2^2(z) & \lambda_3^2(z) \end{pmatrix}$$
(2.4b)

Setting $\tilde{\Psi} = P^{-1}D^{-1}\Psi$, the Lax pair (2.3) becomes

$$\begin{cases}
\tilde{\Psi}_x - q\Lambda(z)\tilde{\Psi} = U\tilde{\Psi}, \\
\tilde{\Psi}_t + (uq\Lambda(z) - \Lambda^{-1}(z))\tilde{\Psi} = V\tilde{\Psi},
\end{cases}$$
(2.5)

where

$$U = \frac{q_x}{3q} \begin{pmatrix} 0 & 1 - \omega^2 & 1 - \omega \\ 1 - \omega & 0 & 1 - \omega^2 \\ 1 - \omega^2 & 1 - \omega & 0 \end{pmatrix}, \tag{2.6a}$$

$$V = -uU + \frac{1}{3z} \left\{ 3\left(\frac{1}{q} - 1\right) \mathbb{I} + \left(q^2 - \frac{1}{q}\right) \begin{pmatrix} 1 & 1 & 1\\ 1 & 1 & 1\\ 1 & 1 & 1 \end{pmatrix} \right\} \begin{pmatrix} \omega^2 & 0 & 0\\ 0 & \omega & 0\\ 0 & 0 & 1 \end{pmatrix}$$
(2.6b)

Define y(x,t) by

$$y(x,t) = \int_{(0,0)}^{(x,t)} q(x',t')(dx' - u(x',t')dt'). \tag{2.7}$$

It is well-defined, because the conservation law

$$q_t + (uq)_x = 0, (2.8)$$

implies that the integral in (2.7) is independent of the path of integration.

Introducing $\tilde{\Phi} = \tilde{\Psi} e^{-y(x,t)\Lambda(z)-t\Lambda^{-1}(z)}$, we have

$$\begin{cases}
\tilde{\Phi}_x - [q\Lambda(z), \tilde{\Phi}] = U\tilde{\Phi}, \\
\tilde{\Phi}_t + [uq\Lambda(z) - \Lambda^{-1}(z), \tilde{\Phi}] = V\tilde{\Phi},
\end{cases}$$
(2.9)

where brackets denote matrix commutator.

We define three contours in (x,t)-domain, see Figure 1,

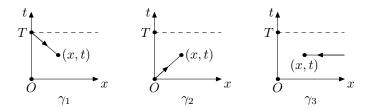


Figure 1: The three contours γ_1, γ_2 and γ_3 in the (x, t)-domain.

And we also denote six sets which decompose the complex z-plane, see Figure 2 In these six sets Ω_n , n = 1, 2, ..., 6, the eigenvalues of $\Lambda(z)$ and

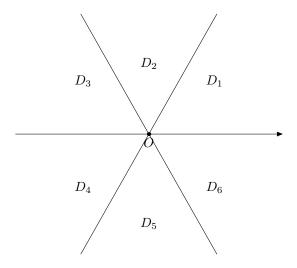


Figure 2: The sets Ω_n , $n=1,2,\ldots,6$, which decompose the complex z-plane.

 $\Lambda^{-1}(z)$ has the following properties,

$$D_{1}: \{\operatorname{Re}\lambda_{1} < \operatorname{Re}\lambda_{2} < \operatorname{Re}\lambda_{3}, \quad \operatorname{Re}\lambda_{1}^{-1} < \operatorname{Re}\lambda_{2}^{-1} < \lambda_{3}^{-1}\},$$

$$D_{2}: \{\operatorname{Re}\lambda_{1} < \operatorname{Re}\lambda_{3} < \operatorname{Re}\lambda_{2}, \quad \operatorname{Re}\lambda_{1}^{-1} < \operatorname{Re}\lambda_{3}^{-1} < \lambda_{2}^{-1}\},$$

$$D_{3}: \{\operatorname{Re}\lambda_{3} < \operatorname{Re}\lambda_{1} < \operatorname{Re}\lambda_{2}, \quad \operatorname{Re}\lambda_{3}^{-1} < \operatorname{Re}\lambda_{1}^{-1} < \lambda_{2}^{-1}\},$$

$$D_{4}: \{\operatorname{Re}\lambda_{3} < \operatorname{Re}\lambda_{2} < \operatorname{Re}\lambda_{1}, \quad \operatorname{Re}\lambda_{3}^{-1} < \operatorname{Re}\lambda_{2}^{-1} < \lambda_{1}^{-1}\},$$

$$D_{5}: \{\operatorname{Re}\lambda_{2} < \operatorname{Re}\lambda_{3} < \operatorname{Re}\lambda_{1}, \quad \operatorname{Re}\lambda_{2}^{-1} < \operatorname{Re}\lambda_{3}^{-1} < \lambda_{1}^{-1}\},$$

$$D_{6}: \{\operatorname{Re}\lambda_{2} < \operatorname{Re}\lambda_{1} < \operatorname{Re}\lambda_{3}, \quad \operatorname{Re}\lambda_{2}^{-1} < \operatorname{Re}\lambda_{1}^{-1} < \lambda_{3}^{-1}\},$$

$$(2.10)$$

The solutions of (2.9) can be constructed as solutions of the Fredholm integral equation

$$(\tilde{\Phi}_{n}(x,t,z))_{ij} = \delta_{ij} + \int_{\gamma_{ij}^{n}} (e^{(y(x,t)-y(x',t'))\hat{\Lambda}(z) + (t-t')\hat{\Lambda}^{-1}(z)} (Udx' + Vdt')(x',t')\tilde{\Phi}(x',t',z))_{ij}, \quad z \in \Omega_{n},$$
(2.11)

where the contours γ_{ij}^n , $n=1,2,\,i,j=1,2,3$ are defined by

$$\gamma_{ij}^{n} = \begin{cases}
\gamma_{1} & if \quad \operatorname{Re}\lambda_{i}(z) < \operatorname{Re}\lambda_{j}(z) \quad and \quad \operatorname{Re}\lambda_{i}^{-1}(z) \ge \operatorname{Re}\lambda_{j}^{-1}(z), \\
\gamma_{2} & if \quad \operatorname{Re}\lambda_{i}(z) < \operatorname{Re}\lambda_{j}(z) \quad and \quad \operatorname{Re}\lambda_{i}^{-1}(z) < \operatorname{Re}\lambda_{j}^{-1}(z), \quad \text{for} \quad z \in \Omega_{n}. \\
\gamma_{3} & if \quad \operatorname{Re}\lambda_{i}(z) \ge \operatorname{Re}\lambda_{j}(z) \quad .
\end{cases}$$
(2.12)

Remark 2.1 For each n = 1, 2, ..., 6, the function $\tilde{\Phi}_n(x, t, z)$ is well-defined by equation (2.11) for $z \in \bar{\Omega}_n$ and (x, t) in the domain (1.7). This is because of the definition of $\{\gamma_j\}_1^3$ and (2.12). The definition of $\{\gamma_j\}_1^3$ implies that

$$\gamma_1: , t - t' \le 0,$$

$$\gamma_2: y(x,t) - y(x',t') \ge 0, t - t' \ge 0,$$

$$\gamma_3: y(x,t) - y(x',t') \le 0.$$
(2.13)

The (ij)th entry of the integral equation (2.18) involves the exponential factor

$$e^{(\lambda_i(z)-\lambda_j(z))(y(x,t)-y(x',t'))+(\lambda_i^{-1}(z)-\lambda_j^{-1}(z))(t-t')}$$

The definition of (2.12) implies that this factor remains bounded for $z \in \Omega_n$ when integrated along the contour γ_{ij}^n . In fact, it is similar to that of $\tilde{\Phi}_{0n}(x,t,z)$, except that since y(x,t) - y(x',t') can take on both signs in the case of γ_1 , the exponential is not necessarily bounded for the integration along γ_1 . However, the matrices $(\gamma^n)_{ij} = \gamma_{ij}^n$ for n = 1, 2, ..., 6 does not involve integration along γ_1 , so we can still conclude that $\tilde{\Phi}_n(x,t,z)$ is well-defined in each Ω_n .

Proposition 2.2 For any fixed point (x,t), $\tilde{\Phi}_n(x,t,z)$ is bounded and analytic as a function of $z \in \Omega_n$, n = 1, 2, ..., 6 away from a possible discrete set of singularities $\{z_j\}$ at which the Fredholm determinant vanishes. Moreover, $\tilde{\Phi}_n(x,t,z)$ admits a bounded and continuous extension to $\bar{\Omega}_n$ and

$$\tilde{\Phi}_n(x,t,z) = \mathbb{I} + O(\frac{1}{z}), \qquad z \to \infty, \quad z \in \Omega_n.$$
 (2.14)

2.1.2 Second Lax pair (well-controlled at z = 0)

Setting $\tilde{\Psi}_0 = P^{-1}\Psi$, then the Lax pair (2.3) becomes

$$\begin{cases}
\tilde{\Psi}_{0x} - \Lambda(z)\tilde{\Psi}_0 = U_0\tilde{\Psi}_0, \\
\tilde{\Psi}_{0t} - \Lambda^{-1}(z)\tilde{\Psi}_0 = V_0\tilde{\Psi}_0,
\end{cases}$$
(2.15)

where

$$U_0 = -\frac{zu_{xx}}{3} \begin{pmatrix} \omega & 0 & 0 \\ 0 & \omega^2 & 0 \\ 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} 1 & 1 & 1 \\ 1 & 1 & 1 \\ 1 & 1 & 1 \end{pmatrix}$$
 (2.16a)

$$V_{0} = \frac{u_{x}}{3} \begin{pmatrix} 0 & 1 - \omega^{2} & 1 - \omega \\ 1 - \omega & 0 & 1 - \omega^{2} \\ 1 - \omega^{2} & 1 - \omega & 0 \end{pmatrix} - zu \begin{pmatrix} \omega & 0 & 0 \\ 0 & \omega^{2} & 0 \\ 0 & 0 & 1 \end{pmatrix} \left\{ \mathbb{I} - \frac{u_{x}x}{3} \begin{pmatrix} 1 & 1 & 1 \\ 1 & 1 & 1 \\ 1 & 1 & 1 \end{pmatrix} \right\}$$

$$(2.16b)$$

Setting $\tilde{\Phi}_0 = \tilde{\Psi}_0 e^{-x\Lambda(z)-t\Lambda^{-1}(z)}$, then the Lax pair of $\tilde{\Phi}_0$ is

$$\begin{cases} \tilde{\Phi}_{0x} - [\Lambda(z), \tilde{\Phi}_0] = U_0 \tilde{\Phi}_0, \\ \tilde{\Phi}_{0t} - [\Lambda^{-1}(z), \tilde{\Phi}_0] = V_0 \tilde{\Phi}_0, \end{cases}$$
 (2.17)

whose solutions can be constructed as solutions of the Fredholm integral equation

$$(\tilde{\Phi}_{0n}(x,t,z))_{ij} = \delta_{ij} + \int_{\gamma_{ij}^n} (e^{(x-x')\hat{\Lambda}(z) + (t-t')\hat{\Lambda}^{-1}(z)} (U_0 dx' + V_0 dt')(x',t') \tilde{\Phi}_0(x',t',z))_{ij}, \quad z \in \Omega_n,$$
(2.18)

where γ_{ij}^n are defined as (2.12).

Remark 2.3 For each n = 1, 2, ..., 6, the function $\tilde{\Phi}_{0n}(x, t, z)$ is well-defined by equation (2.18) for $z \in \bar{\Omega}_n$ and (x, t) in the domain (1.7). This is because of the definition of $\{\gamma_j\}_1^3$ and (2.12). The definition of $\{\gamma_j\}_1^3$ implies that

$$\gamma_1: x - x' \ge 0, t - t' \le 0,$$

$$\gamma_2: x - x' \ge 0, t - t' \ge 0,$$

$$\gamma_3: x - x' \le 0.$$
(2.19)

The (ij)th entry of the integral equation (2.18) involves the exponential factor

$$e^{(\lambda_i(z)-\lambda_j(z))(x-x')+(\lambda_i^{-1}(z)-\lambda_j^{-1}(z))(t-t')}$$

The definition of (2.12) implies that this factor remains bounded for $z \in \Omega_n$ when integrated along the contour γ_{ij}^n .

Proposition 2.4 For any fixed point (x,t), $\tilde{\Phi}_{0n}(x,t,z)$ is bounded and analytic as a function of $z \in \Omega_n$, n = 1, 2, ..., 6 away from a possible discrete set of singularities $\{z_j\}$ at which the Fredholm determinant vanishes. Moreover, $\tilde{\Phi}_{0n}(x,t,z)$ admits a bounded and continuous extension to $\bar{\Omega}_n$ and

$$\tilde{\Phi}_{0n}(x,t,z) = \mathbb{I} + z\tilde{\Phi}_0^{(1)} + z^2\tilde{\Phi}_0^{(2)} + O(z^3), \qquad z \to 0, \quad z \in \Omega_n,$$
 (2.20)

where

$$\tilde{\Phi}_0^{(1)} = -\frac{1}{3}u_x\Gamma, \quad \tilde{\Phi}_0^{(2)} = -\frac{1}{3}u\tilde{\Gamma}, \tag{2.21}$$

$$with \ \Gamma = \left(\begin{array}{ccc} \omega & \omega & \omega \\ \omega^2 & \omega^2 & \omega^2 \\ 1 & 1 & 1 \end{array} \right) \ and \ \tilde{\Gamma} = [\tilde{\Lambda}, \Gamma], \ \tilde{\Lambda} = diag\{\omega, \omega^2, 1\}.$$

2.1.3 Further properties of $ilde{\Phi}_n$ and $ilde{\Phi}_{0n}$

Now, noticing that $\tilde{\Phi}_n$ and $\tilde{\Phi}_{0n}$ are related to the same linear system of PDEs (2.3), tracing back the way that the differential equations for $\tilde{\Phi}_n$ and $\tilde{\Phi}_{0n}$ were derived, leads to the following proposition.

Proposition 2.5 The functions $\tilde{\Phi}_n$ and $\tilde{\Phi}_{0n}$ are related as follows,

$$\tilde{\Phi}_n(x,t,z) = P^{-1}(z)D^{-1}(x,t)P(z)\tilde{\Phi}_{0n}(x,t,z)e^{(x-y(x,t)\Lambda(z))}.$$
(2.22)

where P(z) and D(x,t) defined as (2.4), and y(x,t) is defined as (2.7).

Noticing that it is a straight calculus to show $P^{-1}(z)D^{-1}(x,t)P(z)$ is independent of z.

From (2.14), (2.20) and (2.22) we derive the following expansion of $\tilde{\Phi}_n(x,t,z)$ as $z \to 0$,

$$\tilde{\Phi}_n(x,t,z) = P^{-1}(z)D^{-1}(x,t)P(z)\left(\mathbb{I} + z\left\{-\frac{u_x}{3}\Gamma + (x-y)\tilde{\Lambda}\right\}\right) + z^2\left\{-\frac{u}{3}\tilde{\Gamma} - \frac{u_x}{3}(x-y)\Gamma\tilde{\Gamma} + \frac{(x-y)^2}{2}\tilde{\Lambda}^2\right\} + O(z^3)\right)$$
(2.23)

2.2 Symmetries

Proposition 2.6 $\tilde{\Phi}_n(x,t,z)$ satisfies the symmetry relations:

•
$$\tilde{\Phi}_n(x, t, \omega z) = \mathcal{A}\tilde{\Phi}_n(x, t, z)\mathcal{A}^{-1}$$
, with $\mathcal{A} = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1 & 0 & 0 \end{pmatrix}$,

•
$$\tilde{\Phi}_n(x,t,z) = \mathcal{B}\overline{\tilde{\Phi}_n(x,t,\bar{z})}\mathcal{B}^{-1}$$
, with $\mathcal{B} = \begin{pmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix} = \mathcal{B}^{-1}$,

•
$$\tilde{\Phi}_n(x,t,z) = C\overline{\tilde{\Phi}_n(x,t,\omega^2\bar{z})}C^{-1}$$
, with $C = \begin{pmatrix} 0 & 0 & 1 \\ 0 & 1 & 0 \\ 1 & 0 & 0 \end{pmatrix} = C^{-1}$.

•
$$\tilde{\Phi}_n(x,t,z) = \mathcal{D}\overline{\tilde{\Phi}_n(x,t,\omega\bar{z})}\mathcal{D}^{-1}$$
, with $\mathcal{D} = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix} = \mathcal{D}^{-1}$

The functions $\tilde{\Phi}_{0n}(x,t,z)$ satisfies the same symmetry relations.

3 Sectionally Meromorphic Functions

Let $\{\tilde{\Phi}_n\}_1^6$ and $\{\tilde{\Phi}_{0n}\}_1^6$ denote the eigenfunctions defined in Section 2. We have well control over $\tilde{\Phi}_n$ near $z=\infty$. On the other hand, we have well control over $\tilde{\Phi}_{0n}$ near z=0. Therefore, we will introduce a radius R>0 and formulate a Riemann-Hilbert problem by using the $\tilde{\Phi}_n$ for |z|>R and the $\tilde{\Phi}_{0n}$ for |z|< R.

Define sets $\{D_n\}_1^{12}$ by (see Figure 3)

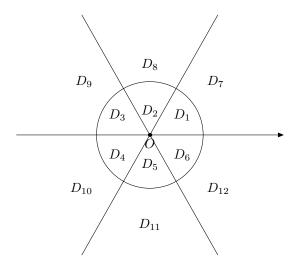


Figure 3: The sets Ω_n , $n=1,2,\ldots,12$, which decompose the complex z-plane.

$$D_n = \Omega_n \cap \{|z| > R\}, \quad n = 1, 2, \dots, 6,$$

$$D_{n+6} = \Omega_n \cap \{|z| < R\}, \quad n = 1, 2, \dots, 6.$$
(3.1)

Since the map $F:(x,t)\to (y,t), y=y(x,t)$ is a bijection from the domain (1.7) onto $F(\Omega)$, we can define functions $\{M_n(y,t,z)\}_1^{12}$ for $(y,t)\in F(\Omega)$ by

$$M_n(y,t,z) = \begin{cases} \tilde{\Phi}_{0n}(z,t,z)e^{(x-y)\Lambda(z)}, & z \in D_n, n = 1, 2, \dots, 6, \\ P^{-1}(z)D(x,t)P(z)\tilde{\Phi}_n(x,t,z), & z \in D_n, n = 7, 8, \dots, 12. \end{cases}$$
(3.2)

The M_n defined in (3.2) are bounded and analytic on the complex z-plane. Reminding that the relation condition between the two eigenfunctions $\tilde{\Phi}_n$ and $\tilde{\Phi}_{0n}$, we need formulating the Riemann-Hilbert problem in terms of the row vectors ν_n defined by

$$\nu_n(y,t,z) = \begin{pmatrix} 1 & 1 & 1 \end{pmatrix} M_n(y,t,z), \quad z \in D_n, \quad n = 1, 2, \dots 12.$$
 (3.3)

Let M and ν denote the sectionally meromorphic functions on the complex z-plane which equal M_n and ν_n respectively for $z \in D_n$.

Lemma 3.1 The function M obeys the symmetries

•
$$M(x, t, \omega z) = \mathcal{A}M(x, t, z)\mathcal{A}^{-1}$$
, with $\mathcal{A} = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1 & 0 & 0 \end{pmatrix}$,

•
$$M(x,t,z) = \mathcal{B}\overline{M(x,t,\bar{z})}\mathcal{B}^{-1}$$
, with $\mathcal{B} = \begin{pmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix} = \mathcal{B}^{-1}$,

•
$$M(x,t,z) = \overline{CM(x,t,\omega^2\bar{z})}C^{-1}$$
, with $C = \begin{pmatrix} 0 & 0 & 1 \\ 0 & 1 & 0 \\ 1 & 0 & 0 \end{pmatrix} = C^{-1}$.

•
$$M(x,t,z) = \mathcal{D}\overline{M(x,t,\omega\bar{z})}\mathcal{D}^{-1}$$
, with $\mathcal{D} = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix} = \mathcal{D}^{-1}$

Proof: This is a consequence of equation (3.2) and the symmetry properties of the $\tilde{\Phi}_n$ and the $\tilde{\Phi}_{0n}$. Noticing that

$$P(z) = P(\omega z)(A), \quad P(z) = \overline{P(\bar{z})}\mathcal{B}.$$
 (3.4)

And the other two symmetries can be obtained by these two symmetries. \Box

3.1 The jump matrices

We define spectral functions $S_n(z)$ by

$$S_n(z) = M_n(0, 0, z), \qquad z \in D_n, \quad n = 1, 2, \dots, 12.$$
 (3.5)

The tracelessness of the matrices $\{U, V\}$ and $\{U_0, V_0\}$ implies that

$$\det S_n(z) = 1, \quad n = 1, 2, \dots, 12. \tag{3.6}$$

The exponential factor $e^{(x-y)\Lambda(z)}$ on the right-hand side of (3.2) has been included because it ensures that the jump matrices introduced in the next proposition depend on x only through the function y(x,t).

Proposition 3.2 For each n = 1, 2, ..., 12, the function ν_n is bounded and analytic in D_n (away from the possibly empty discrete set $\{z_j\}$). Moreover, each ν_n has a continuous and bounded extension to \bar{D}_n . The function ν satisfies the jump conditions

$$\nu_n = \nu_m J_{m,n}, \quad z \in \bar{D}_n \cap \bar{D}_m, \quad n, m = 1, \dots, 12, \quad n \neq m,$$
 (3.7)

where the jump matrix $J_{m,n}(y,t,z)$, $J_{m,n}=J_{n,m}^{-1}$ is defined by

$$J_{m,n} = e^{y\Lambda(z) + t\Lambda^{-1}(z)} (S_m^{-1} S_n) e^{-y\Lambda(z) - t\Lambda^{-1}(z)}, \quad n, m \in \{1, \dots, 12\}.$$
 (3.8)

Proof: The analyticity and boundedness properties of the ν_n follow from the properties of the $\tilde{\Phi}_n$ and the $\tilde{\Phi}_{0n}$ established in Section 2.

From the definition of $\tilde{\Phi}_n$ (2.11) and $\tilde{\Phi}_{0n}$ (2.18), we can deduce that

$$M_n = M_m e^{y\Lambda(z) + t\Lambda^{-1}(z)} J(z) e^{-y\Lambda(z) - t\Lambda^{-1}(z)},$$
 (3.9)

where J(z) is a matrix independent of (x, t).

Evaluation at x = t = 0 yields $J = S_m^{-1}S_n$. Multiplying (3.9) by $\begin{pmatrix} 1 & 1 & 1 \end{pmatrix}$ from the left, we obtain the jump condition (3.7) with $J_{n,m}$ given by (3.8). \square

4 Residue Conditions

If the ν_n have pole singularities at some points $\{z_j\}, z_j \in \mathbb{C}$, the Riemann-Hilbert problem needs to include the residue conditions at these points. We will assume that all z_j lie in the interiors of the sets $\{D_n\}_1^6$, singularities in the interiors of the sets $\{D_n\}_7^{12}$ can be avoided by choosing R large enough. The residue conditions can be found by relating the M_n to another set of solutions of (2.15), denoted by $\{\mu_j\}_1^3$, which are defined by

$$\mu_{j}(x,t,z) = \mathbb{I} + \int_{\gamma_{j}} e^{(x-x')\hat{\Lambda}(z) + (t-t')\Lambda^{-1}(z)} (Udx' + Vdt')(x',t')\mu_{j}(x',t',z), \quad j = 1, 2, 3,$$
(4.1)

where $\{\gamma_j\}_1^3$ are contours shown in Figure 1.

4.1 A matrix factorization problem

Let us define the 3×3 -matrix value spectral functions s(z) and S(z) by

$$\mu_3(x,t,z) = \mu_2(x,t,z)e^{(x\hat{\Lambda}(z)+t\hat{\Lambda}^{-1}(z))}s(z),$$
(4.2a)

$$\mu_1(x,t,z) = \mu_2(x,t,z)e^{(x\hat{\Lambda}(z)+t\hat{\Lambda}^{-1}(z))}S(z).$$
 (4.2b)

Thus,

$$s(z) = \mu_3(0, 0, z), \qquad S(z) = \mu_1(0, 0, z).$$
 (4.3)

Lemma 4.1 Due to the symmetries of M_n , see Lemma 3.1, we just need to calculate S_1 . The S_1 defined in (3.5) can be expressed in terms of the entries of s(z) and S(z) as follows:

$$S_{1} = \begin{pmatrix} s_{11} & 0 & 0 \\ s_{21} & \frac{m_{33}(s)}{s_{11}} & 0 \\ s_{31} & \frac{m_{23}(s)}{s_{11}} & \frac{1}{m_{33}(s)} \end{pmatrix}, \tag{4.4}$$

where m_{ij} denote that the (i, j)-th minor of s.

Proof: Let $\gamma_3^{X_0}$ denote the contour $(X_0,0) \to (x,t)$ in the (x,t)-plane, here $X_0 > 0$ is a constant. We introduce $\mu_3(x,t,z;X_0)$ as the solution of (4.1) with j=3 and with the contour γ_3 replaced by $\gamma_3^{X_0}$. Similarly, we define $M_1(x,t,z;X_0)$ as the solution of (3.2) with γ_3 replaced by $\gamma_3^{X_0}$. We will first derive expression for $S_1(z;X_0) = M_1(0,0,z;X_0)$ in terms of $S_1(z;X_0) = M_2(0,0,z;X_0)$. Then (4.4) will follow by taking the limit $X_0 \to \infty$.

First, from the definition of M_n (3.2), we have the following relations:

$$\begin{cases}
M_1(y,t,z;X_0) = \mu_1(x,t,z)e^{(x\Lambda(z)+t\Lambda^{-1}(z))}R_1(z;X_0)e^{-y\Lambda(z)-t\Lambda^{-1}(z)}, \\
M_1(y,t,z;X_0) = \mu_2(x,t,z)e^{(x\Lambda(z)+t\Lambda^{-1}(z))}S_1(z;X_0)e^{-y\Lambda(z)-t\Lambda^{-1}(z)}, \\
M_1(y,t,z;X_0) = \mu_3(x,t,z;X_0)e^{(x\Lambda(z)+t\Lambda^{-1}(z))}T_1(z;X_0)e^{-y\Lambda(z)-t\Lambda^{-1}(z)}.
\end{cases}$$
(4.5)

Then we get $R_1(z; X_0)$ and $T_1(z; X_0)$ are defined as follows:

$$R_1(z; X_0) = e^{-\hat{\Lambda}^{-1}(z)T} M_1(y(0, T), T, z; X_0) e^{y(0, T)\Lambda(z)}, \tag{4.6a}$$

$$T_1(z; X_0) = e^{-\Lambda(z)X_0} M_n(y(X_0, 0), 0, z; X_0) e^{y(X_0, 0)\Lambda(z)}.$$
 (4.6b)

The relations (4.5) imply that

$$s(z; X_0) = S_1(z; X_0) T_1^{-1}(z; X_0), \qquad S(z) = S_1(z; X_0) R_1^{-1}(z; X_0).$$
 (4.7)

These equations constitute a matrix factorization problem which, given $\{s(z), S(z)\}$ can be solved for the $\{R_1, S_1, T_1\}$. Indeed, the integral equations (3.2) together with the definitions of $\{R_1, S_1, T_1\}$ imply that

$$\begin{cases}
(R_1(z; X_0))_{ij} = 0 & if \quad \gamma_{ij}^1 = \gamma_1, \\
(S_1(z; X_0))_{ij} = 0 & if \quad \gamma_{ij}^1 = \gamma_2, \\
(T_1(z; X_0))_{ij} = 0 & if \quad \gamma_{ij}^1 = \gamma_3.
\end{cases}$$
(4.8)

It follows that (4.7) are 18 scalar equations for 18 unknowns. By computing the explicit solution of this algebraic system, we find that $\{S_1(z; X_0)\}$ are given by the equation obtained from (4.4) by replacing $\{S_1(z), s(z)\}$ with $\{S_1(z; X_0), s(z; X_0)\}$. Taking $X_0 \to \infty$ in this equation, we arrive at (4.4).

4.2 The residue conditions

Since μ_2 is an entire function, it follows that M can only have singularities at the points where the S_1 have singularities. We infer from the explicit formulas (4.4) that the possible singularities of M are as follows:

- $[M]_2$ could have poles in D_1 at the zeros of $s_{11}(z)$;
- $[M]_3$ could have poles in D_1 at the zeros of $m_{33}(z)$.

We denote the above possible zeros by $\{z_j\}_1^N$ and assume they satisfy the following assumption.

Assumption 4.2 We assume that

- $s_{11}(z)$ has n_1 possible simple zeros in D_1 denoted by $\{z_j\}_1^{n_1}$;
- $m_{33}(z)$ has $N-n_1$ possible simple zeros in D_1 denoted by $\{z_j\}_{n_1+1}^N$;

and that none of these zeros coincide. Moreover, we assume that none of these functions have zeros on the boundaries of the D_n 's.

We determine the residue conditions at these zeros in the following:

Proposition 4.3 Let $\{M_1\}$ be the eigenfunctions defined by (3.2) and assume that the set $\{z_j\}_1^N$ of singularities are as the above assumption. Then the following residue conditions hold:

$$Res_{z=z_{j}}[M]_{2} = \frac{m_{33}(s(z_{j}))}{s_{11}(z_{j})s_{21}(z_{j})} e^{\theta_{12}(z_{j})} [M(z_{j})]_{1},$$

$$1 < j \le n_{1}, z_{j} \in D_{1}$$
(4.9a)

$$Res_{z=z_{j}}[M]_{3} = \frac{s_{11}(z_{j})}{m_{33}(\dot{s}(z_{j}))m_{23}(s(z_{j}))} e^{\theta_{23}(z_{j})} [M(z_{j})]_{2},$$

$$n_{1} < j \leq N, z_{j} \in D_{1}$$

$$(4.9b)$$

where $\dot{f} = \frac{df}{dz}$, and θ_{ij} is defined by

$$\theta_{ij}(z) = (\lambda_i - \lambda_j)y + (\lambda_i^{-1} - \lambda_i^{-1})t, \quad i, j = 1, 2, 3.$$
 (4.10)

Proof: We will prove (4.9a). From the relation

$$M_1 = \mu_2 e^{(x\Lambda(z) + t\Lambda^{-1}(z))} S_1(z) e^{-y\Lambda(z) - t\Lambda^{-1}(z)}, \tag{4.11}$$

For i, j = 1, 2, 3, let $\tilde{\theta}_{ij} = (\lambda_i - \lambda_j)x + (\lambda_i^{-1} - \lambda_j^{-1})t$. In view of the expressions for S_1 given in (4.4), the three columns of (4.11) read:

$$[M_1]_1 e^{\lambda_1(y-x)} = [\mu_2]_1 s_{11}(z) + [\mu_2]_2 e^{\tilde{\theta}_{21}} s_{21}(z) + [\mu_2]_3 e^{\tilde{\theta}_{31}} s_{31}(z), \tag{4.12a}$$

$$[M_1]_2 e^{\lambda_1(y-x)} = [\mu_2]_2 \frac{m_{33}(s)}{s_{11}} + [\mu_2]_3 e^{\tilde{\theta}_{32}} \frac{m_{23}(s)}{s_{11}}, \tag{4.12b}$$

$$[M_1]_3 e^{\lambda_1(y-x)} = [\mu_2]_3 \frac{1}{m_{33}(s)}.$$
(4.12c)

In order to prove (4.9a), we suppose that $z_j \in D_1$ is a simple zero of $s_{11}(z)$. Solving (4.12a) and (4.12c) for $[\mu_2]_2$ and $[\mu_2]_3$ and substituting the result in to (4.12b), we find

$$[M_1]_2 = \frac{m_{33}(s)}{s_{11}s_{21}}e^{\theta_{12}}[M_1]_1 - \frac{m_{33}(s)}{s_{21}}e^{\tilde{\theta}_{12} + \lambda_2(x-y)}[\mu_2]_1 + \frac{m_{13}(s)m_{33}(s)}{s_{21}}e^{\theta_{32}}[M_1]_3.$$

Taking the residue of this equation at z_j , we find the condition (4.9a) in the case when $z_j \in D_1$. Similarly, we can get the equation (4.9b). \square

5 The Riemann-Hilbert problem

The sectionally analytic function $\nu(y,t,z)$ defined in section 3 satisfies a Riemann-Hilbert problem which can be formulated in terms of the initial and boundary values of u(x,t). By solving this Riemann-Hilbert problem, the solution of (2.1) can be recovered in parametric form.

Theorem 5.1 Suppose that u(x,t) are a solution of (2.1) on the half-line domain Ω (1.7) with sufficient smoothness and decays as $x \to \infty$. Suppose that the initial and boundary values $\{u_0(x), g_0(t), g_1(t), g_2(t)\}$ defined in (??)

satisfy the assumptions (1.9). Then u(x,t) can be reconstructed from the initial value and boundary values as follows.

Use the initial and boundary data to define $\{\tilde{\Phi}_n(0,0,z)\}_1^6$ and $\{\tilde{\Phi}_{0n}(0,0,z)\}_1^6$ via the integral equations (2.11) and (2.18), respectively. Define spectral functions $S_n(z)$, $n = 1, \ldots, 12$, by

$$S_n(z) = \tilde{\Phi}_{0n-6}(0,0,z), \qquad z \in D_n, \quad n = 1,\dots,6,$$

$$S_n(z) = P^{-1}(z)D(0,0)P(z)\tilde{\Phi}_n(0,0,z), \quad z \in D_n, \quad n = 7,\dots,12,$$
(5.1)

where P(z) and D(x,t) are defined in (2.4). Define the jump matrices $J_{m,n}(y,t,z)$ in terms of the S_n by equation (3.8). Define the spectral s(z) and S(z) by equation (4.2). Assume that the possible zeros $\{z_j\}_1^N$ of the functions $s_{11}(z), m_{33}(s(z))$ are as in assumption 4.2.

Then the solution u(x,t) is given in parametric form by

$$u(x,t) = \tilde{u}(y(x,t),t), \tag{5.2}$$

where

$$x(y,t) = y + \lim_{z \to 0} \frac{1}{z} \left(\nu^{(3)}(y,t,z) - 1 \right),$$
 (5.3a)

$$\tilde{u}(y,t) = \frac{\partial x(y,t)}{\partial t}.$$
 (5.3b)

with $\nu^{(3)}(y,t,z)$ is the third column of row-vector value function $\nu(y,t,z)$ which satisfies the following 3×3 vector Riemann-Hilbert problem:

- $\nu(y,t,z)$ is sectionally meromorphic on the complex z-plane with jumps across the contour $\bar{D}_n \cap \bar{D}_m, n, m = 1, 2, ..., 12$, see Figure 2.
- Across the contour $\bar{D}_n \cap \bar{D}_m$, ν satisfies the jump condition (3.7)
- $\nu(y,t,z) = \begin{pmatrix} 1 & 1 & 1 \end{pmatrix} + O(\frac{1}{z}), \qquad z \to \infty.$
- $\nu^{(2)}$ has simple poles at $z=z_j,\ 1\leq j\leq n_1.$ $\nu^{(3)}$ has simple poles at $z=z_j,\ n_1+1\leq j\leq N.$ The associated residue condition is showed in

the following:

$$Res_{z=z_{j}}\nu^{(2)}(y,t,z) = \frac{m_{33}(s(z_{j}))}{s_{11}(z_{j})s_{21}(z_{j})}e^{\theta_{12}(z_{j})}\nu^{(1)}(y,t,z_{j}),$$

$$1 < j \le n_{1}, z_{j} \in D_{1}$$
(5.4a)

$$Res_{z=z_{j}}\nu^{(3)}(y,t,z) = \frac{s_{11}(z_{j})}{m_{33}(\dot{s}(z_{j}))m_{23}(s(z_{j}))}e^{\theta_{23}(z_{j})}\nu^{(1)}(y,t,z_{j}),$$

$$n_{1} < j \leq N, z_{j} \in D_{1}$$
(5.4b)

where $\dot{f} = \frac{df}{dz}$, and θ_{ij} is defined by

$$\theta_{ij}(z) = (\lambda_i - \lambda_j)y + (\lambda_i^{-1} - \lambda_j^{-1})t, \quad i, j = 1, 2, 3.$$
 (5.5)

• For each zero $z_j \in D_1$, there are five additional points,

$$\omega z_j, \quad \omega^2 z_j, \quad \bar{k}_j, \quad \omega \bar{z}_j, \quad \omega^2 \bar{z}_j,$$

at which ν also has simple poles. The associated residues satisfy the residue conditions obtained from (5.4) and the symmetries of Lemma 3.1.

Proof: The residue conditions (5.4) are obtained by multiplying the conditions in (4.9) by $\begin{pmatrix} 1 & 1 & 1 \end{pmatrix}$ from the left. It remains to prove (5.3). Equation (2.23) implies that

$$\nu(y, t, z) = \begin{pmatrix} \nu_1 & \nu_2 & \nu_3 \end{pmatrix} + O(z^3), \quad as \quad z \to 0,$$
 (5.6)

where

$$\nu_1 = \left(1 + z\omega(x - y) + \frac{1}{2}z^2\omega^2(x - y)^2\right),\tag{5.7a}$$

$$\nu_2 = \left(1 + z\omega^2(x - y) + \frac{1}{2}z^2\omega(x - y)^2\right),\tag{5.7b}$$

$$\nu_3 = \left(1 + z(x - y) + \frac{1}{2}z^2(x - y)^2\right),$$
(5.7c)

Thus, from the equation (5.7c) we can get (5.3). \square

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